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# Properties of Potential Systems in $\mathbf{R}^N$

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## 1. INTRODUCTION

Our purpose is to obtain existence theorems and asymptotic properties as  $|x| \rightarrow \infty$  for potential systems of the type

$$-\Delta u^i + a^i(x) u^i = \lambda f^i(x, \mathbf{u}), \quad x \in \mathbf{R}^N, \quad N \geq 3, \quad i = 1, \dots, M, \quad (1)$$

where  $x = (x^1, \dots, x^N)$ ,  $\mathbf{u} = (u^1, \dots, u^M)$ ,  $\Delta$  denotes the  $N$ -dimensional Laplacian,  $\lambda$  is a positive parameter,  $a^i(x)$  is nonnegative and bounded,  $a^i \in C_{\text{loc}}^\eta(\mathbf{R}^N)$ ,  $f^i \in C_{\text{loc}}^\eta(\mathbf{R}^N \times \mathbf{R}^M, \mathbf{R})$  for some fixed  $\eta \in (0, 1)$ ,  $f^i(x, \mathbf{u}) > 0$  if every  $u^j > 0$ ,  $j = 1, \dots, M$ , and  $f^i$  satisfies additional conditions to be listed in Sections 2 and 3,  $i = 1, \dots, M$ . The definition of a *potential system* (1) is that there exists a differentiable function  $F$  with respect to  $u^1, \dots, u^M$  such that

$$\frac{\partial}{\partial u^i} F(x, u^1, \dots, u^M) = f^i(x, u^1, \dots, u^M), \quad i = 1, \dots, M. \quad (2)$$

In vector-matrix notation, the system (1) is

$$-\Delta \mathbf{u} + A\mathbf{u} = \lambda f(x, \mathbf{u}), \quad x \in \mathbf{R}^N, \quad N \geq 3, \quad (1')$$

where  $A$  denotes the diagonal matrix with diagonal entries  $a^1, \dots, a^M$ . The usual vector length is  $|\mathbf{u}| = [\sum_{i=1}^M |u^i|^2]^{1/2}$  and we also use the notation  $\nabla \mathbf{u} \cdot \nabla \mathbf{v} = \sum_{i=1}^M \nabla u^i \cdot \nabla v^i$ ,  $|\nabla \mathbf{u}|^2 = \nabla \mathbf{u} \cdot \nabla \mathbf{u}$ . Vector inequalities  $\mathbf{u} > \mathbf{v}$  are understood to mean  $u^i > v^i$  for each  $i = 1, \dots, M$ .

The letter  $C$  will be used from time to time as a generic notation for a

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positive constant, whose value is inconsequential and may change from one place to another. The dependence of  $C$  on parameters will be specified only when significant. The open ball in  $\mathbf{R}^N$  of radius  $r$  and centre  $x$  will be denoted by  $B_r(x)$ .

We use the notation  $\|\cdot\|_q, \|\cdot\|_{q,B}$  for the norms in  $L^q(\mathbf{R}^N), L^q(B)$ , respectively, for  $B \subset \mathbf{R}^N$ . The standard Sobolev spaces  $W_0^{1,2}(\mathbf{R}^N)$  and  $D_0^{1,2}(\mathbf{R}^N)$  are the completions of  $C_0^\infty(\mathbf{R}^N)$  in the respective norms

$$\|\phi\|_{1,2} = [\|\nabla\phi\|_2^2 + \|\phi\|_2^2]^{1/2} \quad \text{and} \quad \|\phi\| = \|\nabla\phi\|_2.$$

We also employ Sobolev spaces of vector functions  $\mathbf{u}: \mathbf{R}^M \rightarrow \mathbf{R}$ , defined by

$$E = \{\mathbf{u}: u^i \in W_0^{1,2}(\mathbf{R}^N), i = 1, \dots, M\}$$

and

$$D = \{\mathbf{u}: u^i \in D_0^{1,2}(\mathbf{R}^N), i = 1, \dots, M\}$$

with respective norms

$$\|\mathbf{u}\|_E = [\|\nabla\mathbf{u}\|_2^2 + \|\mathbf{u}\|_2^2]^{1/2} = \left[ \sum_{i=1}^M \|u^i\|_{1,2}^2 \right]^{1/2}$$

and

$$\|\mathbf{u}\|_D = \|\nabla\mathbf{u}\|_2.$$

The existence of a positive solution pair  $\lambda, u$  of (1) will be proved in two cases (i)  $a^i(x)$  is bounded below by a positive constant, and (ii)  $a^i(x) \geq 0$  in  $\mathbf{R}^N, i = 1, \dots, M$ . In both cases  $f^i(x, \mathbf{u})$  is assumed to satisfy a power growth condition in the components of  $\mathbf{u}$ , with exponents  $\beta = (\beta^1, \dots, \beta^M)$  of each power of  $\mathbf{u}$  restricted by  $1 < |\beta| < (N+2)/(N-2)$  (in multi-index notation; see Section 2). Also in both cases  $f^i(x, \mathbf{u}) \rightarrow 0$  as  $|x| \rightarrow \infty$ , uniformly in  $\mathbf{u}$ , but the decay law in case (ii) must be more specific, as described in  $(A_6)$  (see Section 3). The necessity of such conditions for the existence of positive solutions  $\mathbf{u} \in D$  of (1) is indicated in Section 5. None of our results require (1) to be radially symmetric.

In case (i) a direct variational approach is used to construct a positive solution  $\mathbf{u} \in E$  (Theorem 1), but this is not possible in case (ii) since (1) generally has no  $L^2(\mathbf{R}^N)$ -solutions; this is well known if  $M=1, a^1 \equiv 0$ , for example. Instead Theorem 5 yields positive solutions in  $D$  (so in  $L^{2N/(N-2)}(\mathbf{R}^N)$ ) by consideration of a sequence of Eq. (16) $_n, n = 1, 2, \dots$ , of type (i) and letting  $n \rightarrow \infty$ .

A priori asymptotic properties of nonnegative solutions  $\mathbf{u} \in D$  of (1) are proved in Section 4. In particular, conditions on  $f$  are given which guarantee that every such solution satisfies  $u^i(x) = O(|x|^{2-N})$  as  $|x| \rightarrow \infty, i = 1, \dots, M$ . Similar results have been obtained recently by Egnell [12], Li

and Ni [22], and Pucci and Serrin [34] for scalar equations, and sometimes restricted to radially symmetric equations. Of the many recent investigations of problems (1) having  $M = 1$  we mention only [2-6, 8-16, 18-28, 30-32, 36-38, 41]. The few available results for elliptic systems in  $\mathbf{R}^N$  [1, 9, 19, 23, 33] are not of our type here.

## 2. EXISTENCE OF SOLUTIONS IN THE CASE $\mathbf{a}(x) \geq a_* > 0$

The existence of positive solutions  $u$  of (1) will first be proved under the assumption  $a^i(x) \geq a_* > 0$  for all  $x \in \mathbf{R}^N$ ,  $i = 1, \dots, M$ . The assumptions for  $f^i(x, \mathbf{u})$  are listed below. We use the notation

$$f^{ij}(x, \mathbf{u}) = \frac{\partial}{\partial u^j} f^i(x, \mathbf{u}), \quad i, j = 1, \dots, M,$$

and for a multi-index  $\alpha = (\alpha^1, \dots, \alpha^M) \geq 0$ ,

$$\mathbf{u}^\alpha = \prod_{i=1}^M (u^i)^{\alpha^i} \quad \text{if } \mathbf{u} \geq 0, \quad |\alpha| = \sum_{i=1}^M \alpha^i,$$

$$|\mathbf{u}^\alpha| = \prod_{i=1}^M |u^i|^{\alpha^i} \quad \text{if } \mathbf{u} \in \mathbf{R}^M.$$

*Assumptions for  $\mathbf{f}(x, \mathbf{u})$*

(A<sub>1</sub>) There exists a nontrivial potential function  $F$  satisfying (2) such that  $F \in C_{\text{loc}}^{1+\eta}(\mathbf{R}^N \times \mathbf{R}^M, \mathbf{R}_+)$ , where  $\mathbf{R}_+ = [0, \infty)$ ,  $F(x, \mathbf{u})$  is twice continuously differentiable with respect to the components of  $\mathbf{u}$ ,  $F(x, \mathbf{u}) = 0$  if any component  $u^i = 0$ , and  $F(x, -\mathbf{u}) = F(x, \mathbf{u})$  for all  $x, \mathbf{u}$ .

(A<sub>2</sub>) There exists a continuous function  $g: \mathbf{R}^N \rightarrow \mathbf{R}_+$  with  $\lim_{|x| \rightarrow \infty} g(x) = 0$  such that

$$F(x, \mathbf{u}) \leq g(x) \sum_{k=1}^K \mathbf{u}^{\alpha_k}, \quad x \in \mathbf{R}^N, \quad \mathbf{u} \in \mathbf{R}_+^M$$

for some multi-indices  $\alpha_k$  with  $2 < |\alpha_k| < 2N/(N-2)$ ,  $k = 1, \dots, K$ .

(A<sub>3</sub>)  $f^i(x, \mathbf{u}) > 0$  for all  $\mathbf{u} > 0$ ,  $x \in \mathbf{R}^N$  and

$$|f^i(x, \mathbf{u})| \leq Cg(x) \sum_{k=1}^K |\mathbf{u}^{\beta_{ik}}|, \quad x \in \mathbf{R}^N, \quad \mathbf{u} \in \mathbf{R}^M$$

for some constant  $C$  and multi-indices  $\beta_{ik}$  with

$$1 < |\beta_{ik}| < (N+2)/(N-2), \quad i = 1, \dots, M, \quad k = 1, \dots, K.$$

(A<sub>4</sub>)  $|f^{ij}(x, \mathbf{u})| \leq Cg(x) \sum_{k=1}^K |\mathbf{u}^{\gamma_{ijk}}|$ ,  $x \in \mathbf{R}^N$ ,  $\mathbf{u} \in \mathbf{R}^M$  for some constant  $C$  and multi-indices  $\gamma_{ijk}$  satisfying

$$|\gamma_{ijk}| < 4/(N-2), \quad i, j = 1, \dots, M, \quad k = 1, \dots, K.$$

For example, all these conditions hold for systems (1) with potential of the form

$$F(x, \mathbf{u}) = \sum_{k=1}^K g^k(x) |\mathbf{u}^{\alpha_k}|, \quad x \in \mathbf{R}^N, \quad 3 \leq N \leq 5, \quad (3)$$

where each  $g^k \in C^{1+\eta}(\mathbf{R}^N)$ ,  $0 \not\equiv g^k(x) \geq 0$ ,  $\lim_{|x| \rightarrow \infty} g^k(x) = 0$ , either  $\alpha_k^h = 1$  or  $\alpha_k^h > 2$  for each  $h = 1, \dots, M$ ,  $k = 1, \dots, K$ , with  $\alpha_k^h > 2$  for at least one component of each multi-index  $\alpha_k$ , and  $|\alpha_k| < 2N/(N-2)$ ,  $k = 1, \dots, K$ . We note that these restrictions on the multi-indices place serious limitations on the problems (1) which can be considered:  $M \leq 4$  if  $N = 3$ ;  $M \leq 2$  if  $N = 4$  or  $N = 5$ . Such restrictions seem essential for the existence of positive solutions  $\mathbf{u}$  of (1) in  $\mathbf{R}^N$  in view of the necessary conditions in Section 5. If the problem is weakened to seeking nontrivial nonnegative solutions, the restrictions on the multi-indices in (3) can be considerably weakened, e.g., some components of  $\alpha_k$  can be zero.

**THEOREM 1.** *If  $a^i \geq a_* > 0$ ,  $i = 1, \dots, M$  and (A<sub>1</sub>)–(A<sub>4</sub>) hold, the system (1) has a positive solution pair  $\lambda, \mathbf{u}$  with  $\mathbf{u} \in E \cap C_{\text{loc}}^{2+\eta}(\mathbf{R}^N)$ .*

The first step in the proof is to construct a weak solution of (1) by applying the methods of the Calculus of Variations to the constrained minimization problem

$$I^* = \inf\{I(\mathbf{u}) : \mathbf{u} \in E, J(\mathbf{u}) = d > 0\}, \quad (4)$$

where the functionals  $I$  and  $J$  are defined by

$$I(\mathbf{u}) = \frac{1}{2} \int_{\mathbf{R}^N} (|\nabla \mathbf{u}|^2 + A\mathbf{u} \cdot \mathbf{u}) \, dx, \quad \mathbf{u} \in E \quad (5)$$

$$J(\mathbf{u}) = \int_{\mathbf{R}^N} F(x, \mathbf{u}) \, dx, \quad \mathbf{u} \in E. \quad (6)$$

Standard procedure shows that  $I$  is well defined and continuously Fréchet differentiable on  $E$ , with Fréchet derivative  $I'(u)$  given by

$$I'(u) v = \int_{\mathbf{R}^N} (\nabla u \cdot \nabla v + A u \cdot v) dx \quad (7)$$

for all  $u, v \in E$ . An analogue of this for  $J$  will now be proved.

LEMMA 2. *The functional  $J$  is well defined and of class  $C^1(E)$  with Fréchet derivative given by*

$$J'(u) v = \int_{\mathbf{R}^N} \mathbf{f}(x, u) \cdot v dx \quad (8)$$

for all  $u, v \in E$ , where  $\mathbf{f} = (f^1, \dots, f^M)$ .

*Proof.* For  $\alpha = (\alpha^1, \dots, \alpha^M) > 0$  and  $u^h \in L^{|\alpha|}(\mathbf{R}^N)$  for  $h = 1, \dots, M$ , repeated use of Hölder's inequality shows that

$$\left| \int_{\mathbf{R}^N} u^\alpha(x) dx \right| \leq \prod_{h=1}^M \|u^h\|_{|\alpha|}^{\alpha^h}. \quad (9)$$

It follows from (6), (9), and assumption  $(A_2)$  that

$$J(u) \leq C \sum_{k=1}^K \prod_{h=1}^M \|u^h\|_{|\alpha_k|}^{\alpha_k^h}$$

for some positive constant  $C$ . Since  $|\alpha_k| < 2N/(N-2)$  for  $k = 1, \dots, K$ , the Sobolev embedding  $W_0^{1,2}(\mathbf{R}^N) \hookrightarrow L^{|\alpha_k|}(\mathbf{R}^N)$  shows that  $J(u)$  is well defined on  $E$ .

To prove (8) we note from Taylor's formula for  $\phi(t) = F(x, u + tv)$ ,  $u, v \in E$ , and  $(A_4)$  that

$$\begin{aligned} & \left| \frac{J(u + tv) - J(u)}{t} - \int_{\mathbf{R}^N} f(x, u) \cdot v dx \right| \\ &= \left| \frac{1}{t} \int_{\mathbf{R}^N} [F(x, u + tv) - F(x, u) - t \sum_{i=1}^M f^i(x, u) v^i] dx \right| \\ &= \left| \frac{t}{2} \int_{\mathbf{R}^N} \sum_{i,j=1}^M f^{ij}(x, u + \theta tv) v^i v^j dx \right| \\ &\leq C |t| \int_{\mathbf{R}^N} \sum_{i,j=1}^M \sum_{k=1}^K |u + \theta tv|^{y_{ijk}} |v^i| |v^j| dx \end{aligned}$$

for some  $\theta = \theta(x) \in (0, 1)$  and some positive constant  $C$ . We now use (9) with the choice  $\alpha = (\gamma_{ijk}^1, \dots, \gamma_{ijk}^M, 1, 1)$ . By  $(A_4)$ ,  $|\alpha| = |\gamma_{ijk}| + 2 < 2N/(N-2)$  and hence the embedding  $W_0^{1,2}(\mathbf{R}^N) \hookrightarrow L^{|\alpha|}(\mathbf{R}^N)$  shows that  $u^i, v^i \in L^{|\alpha|}(\mathbf{R}^N)$  for  $i = 1, \dots, M$  and  $\|u^h + \theta t v^h\|_{|\alpha|}$  is a bounded function of  $t$  for  $|t| \leq 1$ . Use of (9) for this  $\alpha$  then implies that

$$\begin{aligned} & \left| \frac{J(\mathbf{u} + t\mathbf{v}) - J(\mathbf{u})}{t} - \int_{\mathbf{R}^N} \mathbf{f}(x, \mathbf{u}) \cdot \mathbf{v} \, dx \right| \\ & \leq C |t| \sum_{i,j=1}^M \sum_{k=1}^K \prod_{h=1}^M \|u^h + \theta t v^h\|_{|\alpha|}^{\gamma_{ijk}^h} \|v^i\|_{|\alpha|} \|v^j\|_{|\alpha|} \\ & \rightarrow 0 \quad \text{as} \quad t \rightarrow 0, \end{aligned} \tag{10}$$

thereby establishing (8).

To prove that  $J \in C^1(E)$ , let  $\mathbf{u}, \mathbf{v} \in E$  and let  $\{\mathbf{u}_n\}$  be a sequence in  $E$  with limit  $\mathbf{u}$  in the  $E$ -topology. Let  $\alpha$  be the same multi-index as in (10), so  $|\alpha| < 2N/(N-2)$ . Then the embedding  $W_0^{1,2}(\mathbf{R}^N) \hookrightarrow L^{|\alpha|}(\mathbf{R}^N)$  implies that there exists a constant  $C_0$ , independent of  $n$ , such that

$$\|u_n^j - u^j\|_{|\alpha|} \leq C_0 \|u_n^j - u^j\|_{1,2} \rightarrow 0 \quad \text{as} \quad n \rightarrow \infty$$

for  $j = 1, \dots, M$ . It follows from the mean value theorem,  $(A_4)$ , and (9) by the procedure used for (10) that

$$\begin{aligned} & \left| \int_{\mathbf{R}^N} [\mathbf{f}(x, \mathbf{u}_n) - \mathbf{f}(x, \mathbf{u})] \cdot \mathbf{v} \, dx \right| \\ & \leq C \sum_{i,j=1}^M \sum_{k=1}^K \prod_{h=1}^M \|u^h + \theta(u_n^h - u^h)\|_{|\alpha|}^{\gamma_{ijk}^h} \|v^i\|_{|\alpha|} \|u_n^j - u^j\|_{|\alpha|} \end{aligned}$$

for some constant  $C$  and some  $\theta = \theta(x) \in (0, 1)$ . Therefore,

$$\lim_{n \rightarrow \infty} \sup_{\|v\|_E \leq 1} \left| \int_{\mathbf{R}^N} [\mathbf{f}(x, \mathbf{u}_n) - \mathbf{f}(x, \mathbf{u})] \cdot \mathbf{v} \, dx \right| = 0.$$

This together with (8) proves that  $J \in C^1(E)$  (see [40]).

**LEMMA 3.**  *$J$  is weakly sequentially compact on  $E$ .*

*Proof.* Let  $\{\mathbf{u}_n\}$  be a weakly convergent subsequence of a bounded

sequence in  $E$ , with weak limit  $\mathbf{u} \in E$ . Use of the mean value theorem,  $(A_3)$ , and (9) yield, for some  $\theta = \theta(x) \in (0, 1)$ ,

$$\begin{aligned}
 & |J(\mathbf{u}_n) - J(\mathbf{u})| \\
 & \leq \int_{\mathbf{R}^N} |F(x, \mathbf{u}_n) - F(x, \mathbf{u})| dx \\
 & = \int_{\mathbf{R}^N} \left| \sum_{i=1}^M f^i(x, \mathbf{u} + \theta(\mathbf{u}_n - \mathbf{u}))(u_n^i - u^i) \right| dx \\
 & \leq C \int_{\mathbf{R}^N} g(x) \sum_{i=1}^M \sum_{k=1}^K |[\mathbf{u} + \theta(\mathbf{u}_n - \mathbf{u})]^{\beta_{ik}}| u_n^i - u^i| dx \\
 & \leq C \sum_{i=1}^M \sum_{k=1}^K \prod_{h=1}^M \|u^h + \theta(u_n^h - u^h)\|_{|\beta_{ik}|+1}^{\beta_{ik}^h} \|g(u_n^i - u^i)\|_{|\beta_{ik}|+1}. \quad (11)
 \end{aligned}$$

Since  $|\beta_{ik}| + 1 < 2N/(N-2)$  for each  $i, k$  and  $g(x) \rightarrow 0$  as  $|x| \rightarrow \infty$ , it is known [7, p. 264] that multiplication by  $g$  is a compact mapping from  $W_0^{1,2}(\mathbf{R}^N)$  into  $L^{|\beta_{ik}|+1}(\mathbf{R}^N)$ . Thus  $\{\mathbf{u}_n\}$  has a subsequence  $\{\tilde{\mathbf{u}}_n\}$  such that  $\lim_{n \rightarrow \infty} \|g(\tilde{u}_n^i - u^i)\|_{|\beta_{ik}|+1} = 0$  for  $i = 1, \dots, M$ ,  $k = 1, \dots, K$ . It then follows from (11) that  $J(\tilde{\mathbf{u}}_n) \rightarrow J(\mathbf{u})$  as  $n \rightarrow \infty$ .

*Proof of Theorem 1.* We note by  $(A_1)$  that  $J(\mathbf{u}_0) > 0$  for some  $\mathbf{u}_0 \in E$ , and so we can choose  $d = J(\mathbf{u}_0)$  in (4). Let  $\{\mathbf{u}_n\}$  be a minimizing sequence for (4), i.e.,  $\lim_{n \rightarrow \infty} I(\mathbf{u}_n) = I^*$ ,  $J(\mathbf{u}_n) = d$ , and  $\mathbf{u}_n \geq 0$  without loss of generality by  $(A_1)$ . Since each  $a^i(x)$  is bounded above and below by positive constants,  $I$  is an equivalent norm on  $E$ . Then  $\{\mathbf{u}_n\}$  is bounded in  $E$  and therefore has a weakly convergent subsequence, also denoted by  $\{\mathbf{u}_n\}$ , with weak limit  $\mathbf{u} \in E$ . Lemma 3 shows that this subsequence has a subsequence  $\{\mathbf{u}_n\}$  such that  $\lim_{n \rightarrow \infty} J(\mathbf{u}_n) = J(\mathbf{u}) = d$ . The condition  $J(\mathbf{u}) = 0$  if any  $u^i \equiv 0$ , from  $(A_1)$ , implies that each  $u^i$  is nontrivial and non-negative throughout  $\mathbf{R}^N$ ,  $i = 1, \dots, M$ . It is a standard consequence of the weak convergence of  $\mathbf{u}_n$  to  $\mathbf{u}$  that  $I(\mathbf{u}) = I^*$ . Hence,  $\mathbf{u}$  solves the variational problem (4).

By the Euler-Lagrange method there exists a Lagrange multiplier  $\lambda \in \mathbf{R}$  such that  $I'(\mathbf{u}) = \lambda J'(\mathbf{u})$  in  $E^*$ , i.e., by (7) and (8),

$$\int_{\mathbf{R}^N} (\nabla \mathbf{u} \cdot \nabla \mathbf{v} + A \mathbf{u} \cdot \mathbf{v}) dx = \lambda \int_{\mathbf{R}^N} \mathbf{f}(x, \mathbf{u}) \cdot \mathbf{v} dx \quad (12)$$

for all  $\mathbf{v} \in E$ . With the choice  $v^i = \delta_{ij} u^j$ , (12) becomes

$$\int_{\mathbf{R}^N} (|\nabla u^i|^2 + a^i(u^i)^2) dx = \lambda \int_{\mathbf{R}^N} f^i(x, \mathbf{u}) u^i dx, \quad (13)$$

for  $i = 1, \dots, M$ . Since  $0 \neq u^i \geq 0$  in  $\mathbf{R}^N$  and  $f^i(x, \mathbf{u}) > 0$  if  $\mathbf{u} > 0$ , it follows that  $\lambda > 0$ . Equation (12) means that  $\mathbf{u}$  is a weak solution of the system (1). By the regularity hypotheses on  $\mathbf{a}$  and  $\mathbf{f}$ , a standard bootstrap argument shows that each  $u^i \in C_{\text{loc}}^{2+\eta}(\mathbf{R}^N)$ . The strong maximum principle for  $-\Delta u^i + a^i u^i \geq 0$  in  $\mathbf{R}^N$  implies that each  $u^i > 0$  in  $\mathbf{R}^N$ .

**COROLLARY 4.** *The solution  $\mathbf{u}$  in Theorem 1 has the asymptotic behavior  $\lim_{|x| \rightarrow \infty} \mathbf{u}(x) = 0$  and  $\lim_{|x| \rightarrow \infty} |(\nabla \mathbf{u})(x)| = 0$ .*

*Proof.* Let  $U = \sum_{i=1}^M u^i$ . Then  $U$  satisfies the equation  $-\Delta U = \phi(x)$ , where by (A<sub>3</sub>)

$$\begin{aligned} |\phi(x)| &= \left| \sum_{i=1}^M [\lambda f^i(x, \mathbf{u}(x)) - a^i(x) u^i(x)] \right| \\ &\leq C g(x) \sum_{i,k} [U(x)]^{|\beta_{ik}|} + a^* U(x). \end{aligned}$$

Since  $|\beta_{ik}| < (N+2)/(N-2)$  a bootstrap argument as in [30, Lemma 4.4] yields the estimate

$$\|U\|_{C^{1+\eta}(\overline{B_1(x)})} \leq C \|U\|_{1,2,B_2(x)}$$

for a constant  $C$  independent of  $x \in \mathbf{R}^N$ , where  $B_r(x)$  denotes the ball in  $\mathbf{R}^N$  of radius  $r$  and centre  $x$ . Since each  $u^i > 0$ , this implies the conclusion of Corollary 4.

### 3. EXISTENCE IN THE CASE $\mathbf{a} \geq 0$

To obtain an analogue of Theorem 1 when the components of  $\mathbf{a}(x)$  are not bounded below by a positive constant, e.g., when some of the components are identically zero, we need to sharpen assumptions (A<sub>2</sub>) and (A<sub>3</sub>) for  $\mathbf{u} > 0$  by adjoining the following conditions:

(A<sub>5</sub>) There exists a constant  $C > 0$  such that  $u^i f^i(x, \mathbf{u}) \geq CF(x, \mathbf{u})$  for all  $\mathbf{u} \geq 0$ ,  $x \in \mathbf{R}^N$ ,  $i = 1, \dots, M$ .

(A<sub>6</sub>) These exist continuous functions  $g^k: \mathbf{R}^N \rightarrow \mathbf{R}_+$  with  $g^k(x) = 0(|x|^{-b_k})$  as  $|x| \rightarrow \infty$ ,  $b_k > 0$ ,  $k = 1, \dots, K$ , such that

$$0 \leq f^i(x, \mathbf{u}) \leq \sum_{k=1}^K g^k(x) \mathbf{u}^{\beta_k}, \quad x \in \mathbf{R}^N, \quad \mathbf{u} \in \mathbf{R}_+^M$$



for multi-indices  $\beta_{ik}$  satisfying

$$\begin{cases} 1 < |\beta_{ik}| < \frac{N+2}{N-2} & \text{if } b_k \geq 2 \\ \frac{N+2-2b_k}{N-2} < |\beta_{ik}| < \frac{N+2}{N-2} & \text{if } 0 < b_k < 2 \end{cases}$$

$i = 1, \dots, M, k = 1, \dots, K$ .

Conditions  $(A_5)$  and  $(A_6)$  together imply the existence of multi-indices  $\alpha_k$  such that

$$F(x, \mathbf{u}) \leq \frac{1}{C} \sum_{k=1}^K g^k(x) \mathbf{u}^{\alpha_k}, \quad x \in \mathbf{R}^N, \quad \mathbf{u} \in \mathbf{R}_+^M, \quad (14)$$

where

$$\begin{aligned} 2 < |\alpha_k| < \frac{2N}{N-2} & \quad \text{if } b_k \geq 2 \\ \frac{2N-2b_k}{N-2} < |\alpha_k| < \frac{2N}{N-2} & \quad \text{if } 0 < b_k < 2 \end{aligned} \quad (15)$$

for  $k = 1, \dots, K$ .

The necessity of the conditions on  $|\beta_{ik}|$  in  $(A_6)$  will be indicated in Section 5. This, of course, is known in the scalar case  $M = 1$  [11, 12, 16, 18, 21, 22, 27, 32].

**THEOREM 5.** *If  $a^i \geq 0$  in  $\mathbf{R}^N$  and  $(A_1)$ – $(A_6)$  hold, the system (1) has a positive solution pair  $\lambda, \mathbf{u}$  with  $\mathbf{u} \in D \cap C_{\text{loc}}^{2+\eta}(\mathbf{R}^N)$  such that  $\mathbf{u}(x) = 0(|x|^{(2-N)/2})$  as  $|x| \rightarrow \infty$ . If  $a^i$  is identically zero for  $i = 1, \dots, M$ , this asymptotic estimate can be replaced by  $\mathbf{u}(x) = 0(|x|^{2-N})$  as  $|x| \rightarrow \infty$ .*

The idea of the proof is to apply Theorem 1 to the sequence of systems

$$-\Delta u^i + \left[ a^i(x) + \frac{1}{n} \right] u^i = \lambda f^i(x, \mathbf{u}), \quad x \in \mathbf{R}^N, \quad n = 1, 2, \dots, \quad (16)_n$$

of type (1). We are going to prove, for a positive solution pair  $\lambda_n, \mathbf{u}_n$  of  $(16)_n$ , that  $\{\lambda_n\}, \{\mathbf{u}_n\}$  have subsequences that converge to a positive solution pair  $\lambda, \mathbf{u}$  of (1) as  $n \rightarrow \infty$ , with  $\mathbf{u} \in D \cap C_{\text{loc}}^{2+\eta}(\mathbf{R}^N)$ . If each  $a^i \equiv 0$ , the estimate  $\mathbf{u}(x) = 0(|x|^{2-N})$  as  $|x| \rightarrow \infty$  will then follow from Theorem 9 (in Section 4). The analogue of (5) for Eq.  $(16)_n$  is

$$I_n(\mathbf{u}) = \frac{1}{2} \int_{\mathbf{R}^N} \left[ |\nabla \mathbf{u}|^2 + A \mathbf{u} \cdot \mathbf{u} + \frac{1}{n} |\mathbf{u}|^2 \right] dx, \quad \mathbf{u} \in E, \quad n = 1, 2, \dots \quad (17)$$

LEMMA 6. *Under the assumptions of Theorem 5, for every  $n = 1, 2, \dots$ , the system  $(16)_n$  has a positive solution pair  $\lambda_n, \mathbf{u}_n$ , with  $\mathbf{u}_n \in E \cap C_{\text{loc}}^{2+\eta}(\mathbf{R}^N)$ , such that the sequences  $\{\lambda_n\}$  and  $\{\|\nabla \mathbf{u}_n\|_2\}$  are bounded.*

*Proof.* Theorem 1 guarantees the existence of such a pair  $\lambda_n, \mathbf{u}_n$ , where  $\mathbf{u}_n$  solves the variational problem

$$I_n^* = \inf\{I_n(\mathbf{u}) : \mathbf{u} \in E, J(\mathbf{u}) = d > 0\}. \quad (18)$$

It follows from (17) and (18) that

$$I_n(\mathbf{u}_n) = I_n^* \leq \inf\{I_1(\mathbf{u}) : \mathbf{u} \in E, J(\mathbf{u}) = d\} = I_1(\mathbf{u}_1),$$

implying the boundedness of  $\{\|\nabla \mathbf{u}_n\|_2\}$ . The analogue of (13) for  $\lambda_n, \mathbf{u}_n$  is

$$\begin{aligned} \int_{\mathbf{R}^N} \left[ |\nabla u_n^i|^2 + \left[ a^i + \frac{1}{n} \right] (u_n^i)^2 \right] dx &= \lambda_n \int_{\mathbf{R}^N} f^i(x, \mathbf{u}_n) u_n^i dx, \\ i &= 1, \dots, M, \quad n = 1, 2, \dots \end{aligned} \quad (19)$$

As a consequence of (17)–(19) and (A<sub>5</sub>) we obtain

$$\begin{aligned} 2I_1(\mathbf{u}_1) &\geq 2I_n(\mathbf{u}_1) \geq 2I_n(\mathbf{u}_n) > \lambda_n \int_{\mathbf{R}^N} f^i(x, \mathbf{u}_n) u_n^i dx \\ &\geq C\lambda_n \int_{\mathbf{R}^N} F(x, \mathbf{u}_n) dx = C\lambda_n J(\mathbf{u}_n) = C d\lambda_n, \end{aligned}$$

showing that  $\{\lambda_n\}$  is bounded.

LEMMA 7. *The sequence  $\{\mathbf{u}_n(x)\}$  in Lemma 6 is uniformly bounded in  $\mathbf{R}^N$ . Furthermore, there exists a constant  $C$ , independent of  $n$ , such that*

$$u_n^i(x) \leq C |x|^{(2-N)/2}, \quad |x| \geq 1, \quad i = 1, \dots, M, \quad n = 1, 2, \dots \quad (20)$$

*Proof.* From  $(16)_n$ , (A<sub>6</sub>), and the boundedness of  $\{\lambda_n\}$  (by Lemma 6), the function defined by  $U_n = \sum_{i=1}^M u_n^i$  is a subsolution of the linear equation

$$-\Delta U = Ch_n(x) U, \quad x \in \mathbf{R}^N, \quad n = 1, 2, \dots \quad (21)$$

for some positive constant  $C$ , independent of  $n$ , where

$$h_n(x) = \sum_{i=1}^M \sum_{k=1}^K g^k(x) [U_n(x)]^{|\beta_{ik}| - 1}. \quad (22)$$

For  $|x| = 2R \geq 1$  and a fixed constant  $s > N/2$ , we define

$$\phi(R, n) = R^{2 - N/s} \|h_n\|_{s, B_R(x)}.$$

Trudinger's a priori estimate [39, Theorem 5.1] for subsolutions of (21) implies that

$$\sup_{y \in B_{R/2}(x)} U_n(y) \leq C_0 R^{(2 - N)/2} \|U_n\|_{2N/(N-2), B_R(x)} \quad (23)$$

for a constant  $C_0$  depending only on  $N$ ,  $s$ , and  $\phi(R, n)$ . Since  $\|\nabla U_n\|_2$  is bounded by Lemma 6, also  $\|U_n\|_{2N/(N-2)}$  is bounded by a Sobolev embedding theorem. It will be shown below that  $\phi(R, n)$  is uniformly bounded for  $2R \geq 1$ . We can then conclude from (23) that there exists a constant  $C$ , independent of  $n$ , such that

$$U_n(x) \leq CR^{(2 - N)/2} \quad \text{for } |x| = 2R \geq 1,$$

implying (20).

For  $R = 1$ , for example, (23) also implies that  $U_n(x)$  is locally uniformly bounded, and hence each  $u_n^i(x)$  is uniformly bounded in  $\mathbf{R}^N$  by (20),  $i = 1, \dots, M$ .

To prove that  $\phi(R, n)$  is uniformly bounded for  $2R \geq 1$ , we choose  $s$  satisfying

$$\frac{N}{2} < s < \frac{2N}{(N-2)(\beta-1)}, \quad \beta = \max_{i,k} |\beta_{ik}|, \quad (24)$$

possible since  $(N-2)(\beta-1) < 4$  by  $(A_6)$ . Let Hölder exponents  $p_{ik}$ ,  $q_{ik}$  be defined by

$$p_{ik} = \frac{2N}{s(N-2)(|\beta_{ik}| - 1)}, \quad q_{ik} = \frac{2N}{2N - s(N-2)(|\beta_{ik}| - 1)}. \quad (25)$$

Then  $p_{ik} > 1$ ,  $q_{ik} > 1$ , and  $p_{ik}^{-1} + q_{ik}^{-1} = 1$ . Hölder's inequality applied to (22) yields

$$\|h_n\|_{s, B_R(x)}^s \leq (MK)^{s-1} \sum_{i,k} \|(g^k)^s\|_{q_{ik}, B_R(x)} \|U_n^{s(|\beta_{ik}| - 1)}\|_{p_{ik}, B_R(x)}. \quad (26)$$

Note from (25) that  $s(|\beta_{ik}| - 1) p_{ik} = 2N/(N-2)$ . Since  $\|U_n\|_{2N/(N-2)}$  is uniformly bounded, it then follows from  $(A_6)$  and (26) that there exists a constant  $C$ , independent of  $n$ , such that

$$[\phi(R, n)]^s = R^{2s - N} \|h_n\|_{s, B_R(x)}^s \leq C \sum_{i,k} R^{\sigma_{ik}},$$

where

$$\sigma_{ik} = 2s - N - sb_k + \frac{N}{q_{ik}} = \frac{s}{2} [4 - 2b_k - (N-2)(|\beta_{ik}| - 1)].$$

Since  $\sigma_{ik} < 0$  for all  $i, k$  by  $(A_6)$ ,  $\phi(R, n)$  is uniformly bounded for  $2R \geq 1$ . This completes the proof of Lemma 7.

*Proof of Theorem 5.* Let  $\{\lambda_n\}$ ,  $\{\mathbf{u}_n\}$  be the sequences in Lemma 6. Since each  $\{u_n^i\}$ ,  $i = 1, \dots, M$  is uniformly bounded in  $\mathbf{R}^N$  by Lemma 7, standard elliptic regularity theory applied to the  $i$ th component Eq. (1) shows that  $\{u_n(x)\}$  has a subsequence which converges locally uniformly in  $C^2(\mathbf{R}^N)$  to a function  $u^i \in C_{\text{loc}}^2(\mathbf{R}^N)$  satisfying

$$u^i(x) \leq C |x|^{(2-N)/2}, \quad |x| \geq 1, \quad i = 1, \dots, M. \quad (27)$$

In view of Lemma 6,  $\{\lambda_n\}$  has a convergent subsequence with limit  $\lambda \geq 0$ . We can therefore let  $n \rightarrow \infty$  in  $(16)_n$  to conclude that  $\mathbf{u}$  is a nonnegative solution vector of the system (1) whose components satisfy (27).

To prove that each  $u^i(x)$  is not identically zero, we first note from (15) that  $(N-2)|\alpha_k| = 2N - 2b_k + e_k$  for some  $e_k > 0$ ,  $k = 1, \dots, K$ . Exponents for Hölder's inequality will now be selected to be

$$p_k = \frac{2N + e_k}{(N-2)|\alpha_k|}, \quad q_k = \frac{2N + e_k}{2N + e_k - (N-2)|\alpha_k|}. \quad (28)$$

Then  $b_k q_k > N$ ,  $p_k > 1$ ,  $q_k > 1$ , and  $p_k^{-1} + q_k^{-1} = 1$ . Use of (6), (14), (18), and  $(A_6)$  shows that there exists a positive constant  $C$ , independent of  $n$ , such that

$$0 < d = J(\mathbf{u}_n) \leq C \int_{\mathbf{R}^N} \sum_{k=1}^K [1 + |x|^{b_k}]^{-1} [\mathbf{u}_n(x)]^{\alpha_k} dx. \quad (29)$$

Let  $\delta_k$  be a multi-index with the properties  $0 < \delta_k^i < \alpha_k^i$  for  $i = 1, \dots, M$  and  $|\delta_k| = e_k/(N-2)p_k$ ,  $k = 1, \dots, K$ . This is possible since

$$|\alpha_k| - |\delta_k| = \frac{2N}{(N-2)p_k} > 0, \quad k = 1, \dots, K.$$

We define functions  $\phi_{nk}$  and  $\psi_{nk}$  in  $\mathbf{R}^N$  by

$$\phi_{nk}(x) = [\mathbf{u}_n(x)]^{\alpha_k - \delta_k}, \quad \psi_{nk}(x) = [1 + |x|^{b_k}]^{-1} [\mathbf{u}_n(x)]^{\delta_k},$$

$k = 1, \dots, K$ ;  $n = 1, 2, \dots$ . Hölder's inequality applied to (29) yields

$$0 < d \leq \tilde{C} \sum_k \|\phi_{nk}\|_{p_k} \|\psi_{nk}\|_{q_k} \quad (30)$$

for some constant  $\tilde{C}$ , independent of  $n$ . By (9),

$$\|\phi_{nk}\|_{p_k}^{p_k} \leq \prod_{h=1}^M \|u_n^h\|_{p_k}^{p_k(\alpha_k^h - \delta_k^h)},$$

and hence it follows from  $p_k |\alpha_k - \delta_k| = 2N/(N-2)$  that  $\{\|\phi_{nk}\|_{p_k}\}$  is a bounded sequence,  $k = 1, \dots, K$ . Since  $\{u_n(x)\}$  is uniformly bounded by Lemma 7 and  $b_k q_k > N$ , also  $\{\|\psi_{nk}\|_{q_k}\}$  is bounded. From the pointwise convergence of  $\psi_{nk}(x)$  to the function

$$\psi_k(x) = [1 + |x|^{b_k}]^{-1} [u(x)]^{\delta_k}$$

as  $n \rightarrow \infty$ , it follows that  $\psi_k \in L^1(\mathbf{R}^N)$ . Then the dominated convergence theorem applied to (30) gives, for some positive constant  $C$ ,

$$0 < d \leq C \sum_{k=1}^K \|\psi_k\|_{q_k},$$

which is impossible if  $u^i(x)$  is identically zero for any integer  $i = 1, \dots, M$ .

Since  $(\nabla u_n^i)(x) \rightarrow (\nabla u^i)(x)$  pointwise in  $\mathbf{R}^N$  and  $\{\|\nabla u_n^i\|_2\}$  is bounded, Fatou's lemma implies that  $u^i \in D_0^{1,2}(\mathbf{R}^N)$  for  $i = 1, \dots, M$ , i.e.,  $\mathbf{u} \in D$ . If  $\lambda = 0$ ,  $u^i(x)$  would be a nontrivial nonnegative solution of  $\Delta u^i = 0$  in  $\mathbf{R}^N$  with uniform limit zero as  $|x| \rightarrow \infty$  by (27), contrary to Liouville's theorem, and hence  $\lambda > 0$ . Then the strong maximum principle for  $\Delta u^i \leq 0$  shows that  $u^i(x) > 0$  throughout  $\mathbf{R}^N$ ,  $i = 1, \dots, M$ . If each  $a^i \equiv 0$ , Theorem 9 (in Section 4) implies that each  $u^i(x) = 0(|x|^{2-N})$  as  $|x| \rightarrow \infty$ . This completes the proof of Theorem 5.

#### 4. A PRIORI ASYMPTOTIC ESTIMATES

This section concerns the asymptotic behavior of any nontrivial nonnegative solution  $\mathbf{u}$  of (1) such that each component  $u^i \in D_0^{1,2}(\mathbf{R}^N)$ . Similar estimates for the case  $M = 1$  were obtained recently by Egnell [12] via a different procedure.

**THEOREM 8.** *Suppose  $f^i$  satisfies  $(A_3)$  as well as the conditions stated below (1),  $i = 1, \dots, M$ . If  $\mathbf{u} \in D$  is a nonnegative solution of (1), then  $\mathbf{u}$  is bounded,  $\mathbf{u} \in C_{\text{loc}}^2(\mathbf{R}^N)$ , and*

$$\lim_{|x| \rightarrow \infty} \mathbf{u}(x) = \lim_{|x| \rightarrow \infty} |\nabla \mathbf{u}(x)| = 0. \quad (31)$$

*Proof.* The function  $U = \sum_{i=1}^M u^i$  (assumed to be nontrivial) satisfies the linear equation  $-\Delta U = H(x) U$  in  $\mathbf{R}^N$ , where by  $(A_3)$

$$H(x) \leq C \sum_{i=1}^M \sum_{k=1}^K [U(x)]^{|\beta_{ik}| - 1}, \quad (32)$$

similarly to (22). Since  $U \in L^{2N/(N-2)}(\mathbf{R}^N)$  from the embedding  $D_0^{1,2}(\mathbf{R}^N) \hookrightarrow L^{2N/(N-2)}(\mathbf{R}^N)$  and since  $|\beta_{ik}| - 1 < 4/(N-2)$ , Hölder's inequality with exponents

$$p_k = \frac{4}{(N-2)(|\beta_{ik}| - 1)}, \quad q_k = \frac{4}{4 - (N-2)(|\beta_{ik}| - 1)}$$

shows that  $\|H\|_{N/2, B_r(x)} \rightarrow 0$  as  $r \rightarrow 0$  uniformly in  $\mathbf{R}^N$ . Results of Brézis and Kato [8] then imply that  $U \in L^q(\mathbf{R}^N)$  for all  $q \geq 2N/(N-2)$ , and hence the norms

$$\|U\|_{q, B_2(x)}, \|G\|_{q/\beta, B_2(x)} \quad (33)$$

are bounded functions of  $x \in \mathbf{R}^N$  and have limits zero as  $|x| \rightarrow \infty$ , where  $G = HU$  and  $\beta = \min_{i,k} |\beta_{ik}|$ . Standard a priori interior estimates [17, Chap. 8] for  $-\Delta U = G$  show, if  $q$  is sufficiently large in (33), that  $U(x)$  is bounded in  $\mathbf{R}^N$  and  $U \in C_{\text{loc}}^2(\mathbf{R}^N)$ . Then  $\|U\|_{2, B_2(x)}$ , as well as the norms (33), has limit zero as  $|x| \rightarrow \infty$  (by Hölder's inequality), and consequently interior Hölder estimates for  $-\Delta U = G$  [17, Theorems 8.24 and 8.32] imply that

$$\lim_{x \rightarrow \infty} \|U\|_{C^{1+\eta}(\overline{B_1(x)})} = 0.$$

Since each  $u^i \geq 0$  this proves (31).

The next theorem concerns the a priori asymptotic behavior at  $\infty$  of solutions  $\mathbf{u}$  of systems of type

$$-\Delta u^i = f^i(x, \mathbf{u}), \quad x \in \Omega, \quad i = 1, \dots, M \quad (34)$$

in exterior domains  $\Omega \subset \mathbf{R}^N$ , which for our purpose can be taken as  $\Omega = \{x \in \mathbf{R}^N : |x| > R\}$ ,  $R > 0$ . We assume that each  $f^i \in C_{\text{loc}}^\eta(\Omega \times \mathbf{R}^M, \mathbf{R}_+)$ ,  $0 < \eta < 1$ , and that  $f^i$  satisfies the growth condition  $(A_6)$  (see Section 3) in  $\Omega$ , i.e.,  $\mathbf{R}^N$  in  $(A_6)$  is replaced by  $\Omega$ . The notation  $\mathbf{u} \in D(\Omega)$  is defined to mean that  $h\mathbf{u} \in D$  for some smooth nonnegative function  $h$  with  $h(x) = 1$  for all sufficiently large  $x$ .

**THEOREM 9.** *Under the above conditions on  $\mathbf{f}$ , every positive solution  $\mathbf{u} \in D(\Omega)$  of (34) is minimal, i.e.,  $|x|^{N-2} u^i(x)$  is bounded above and below by positive constants in  $\Omega$ ,  $i = 1, \dots, M$ .*

*Proof.* Kelvin's transformation

$$y = x/|x|^2, \quad v^i(y) = |x|^{N-2} u^i(x), \quad i = 1, \dots, M \quad (35)$$

maps (34) into

$$-\Delta v^i = \frac{1}{|y|^{N+2}} f^i \left( \frac{y}{|y|^2}, |y|^{N-2} \mathbf{v}(y) \right), \quad y \in \Omega', \quad i = 1, \dots, M,$$

where  $\Omega' = \{y \in \mathbf{R}^N : 0 < |y| < 1/R\}$ . Then  $v^i$  can be regarded as a solution of the linear equation

$$-\Delta v^i = H^i(y) v^i, \quad y \in \Omega', \quad (36)$$

where

$$H^i(y) = |y|^{-N-2} [v^i(y)]^{-1} f^i \left( \frac{y}{|y|^2}, |y|^{N-2} \mathbf{v}(y) \right).$$

It follows from assumption (A<sub>6</sub>) that

$$H^i(y) \leq \sum_{k=1}^K H^{ik}(y), \quad y \in \Omega',$$

where

$$H^{ik}(y) = |y|^{\rho_{ik}} g^k \left( \frac{y}{|y|^2} \right) (v^1(y))^{\beta_{ik}^1} \dots (v^i(y))^{\beta_{ik}^i-1} \dots (v^M(y))^{\beta_{ik}^M} \quad (37)$$

and  $\rho_{ik} = |\beta_{ik}|(N-2) - N - 2$ ,  $i = 1, \dots, M$ ,  $k = 1, \dots, K$ . Let  $V(y) = \sum_{i=1}^M v^i(y)$ . Then (37) and (A<sub>6</sub>) show that a positive constant  $C$  exists such that

$$H^{ik}(y) \leq C |y|^{\rho_{ik} + b_k} [V(y)]^{|\beta_{ik}|-1}, \quad y \in \Omega'. \quad (38)$$

If we can show that  $H^{ik} \in L^s(\Omega')$  for some  $s > N/2$ ,  $i = 1, \dots, M$ ,  $k = 1, \dots, K$ , then  $H^i \in L^s(\Omega')$ , and Serrin's theorem on isolated singularities [35, p. 220] applied to (36) near  $y = 0$  shows that either  $v^i(y)$  or  $|y|^{N-2} v^i(y)$  is bounded above and below by positive constants in a deleted neighbourhood of  $y = 0$ . Thus either  $|x|^{N-2} u^i(x)$  or  $u^i(x)$  is bounded above and below by positive constants in  $\Omega$ . However, the second possibility is precluded by the assumption that  $\mathbf{u} \in D(\Omega)$ , proving the minimality of  $\mathbf{u}$ .

To show that  $H^{ik} \in L^s(\Omega')$  for  $s = N/(2 - \varepsilon)$ , we fix  $\varepsilon > 0$  satisfying

$$4 - 2b_k + 2\varepsilon < (N-2)(|\beta_{ik}| - 1) < 4 - 2\varepsilon,$$

which is possible by (A<sub>6</sub>). This is equivalent to

$$8 - 2b_k - 2N/s < (N-2)(|\beta_{ik}| - 1) < 2N/s. \quad (39)$$

Hölder's inequality with exponents

$$p_{ik} = \frac{2N}{s(N-2)(|\beta_{ik}| - 1)}, \quad q_{ik} = \frac{2N}{2N - s(N-2)(|\beta_{ik}| - 1)}$$

applied to (38) gives

$$\|H^{ik}\|_{s, \Omega'}^s \leq \| |y|^{s(\rho_{ik} + b_k)} \|_{q_{ik}, \Omega'} \|V\|_{2N/(N-2), \Omega'}^{s(|\beta_{ik}| - 1)}. \quad (40)$$

It is well known and easily verified that the assumption  $\mathbf{u} \in D(\Omega)$  implies that  $h'v \in D$  for some smooth nonnegative function  $h'$  with  $h'(y) = 1$  in a neighborhood of the origin, from which  $V \in L^{2N/(N-2)}(\Omega')$  by Sobolev embedding. Since the left inequality (39) is equivalent to

$$s(\rho_{ik} + b_k) q_{ik} > -N$$

by a routine computation, the conclusion  $H^{ik} \in L^s(\Omega')$  is a consequence of (40). This completes the proof of Theorem 9.

*Remarks.* If  $a^i$  is identically zero for each  $i = 1, \dots, M$ , Theorem 9 can be applied to the solution  $\mathbf{u}(x)$  of (1) obtained in Theorem 5, as already mentioned in Section 3. Elliptic regularity theory for (36) enables us to prove that the function  $v^i(y)$  in (35) actually has a positive limit as  $|y| \rightarrow 0$ , from which  $|x|^{N-2} u^i(x)$  has a positive finite limit as  $|x| \rightarrow \infty$  for  $i = 1, \dots, M$ ; see, e.g., [12].

## 5. NECESSARY CONDITIONS FOR POSITIVE SOLUTIONS

Application of the divergence theorem to  $(x \cdot \nabla u^i) \nabla u^i$  in a ball of radius  $R$ , summing over  $i$ , and then letting  $R \rightarrow \infty$ , we find that every solution  $\mathbf{u} \in E$  of (1) satisfies the Pohožaev-type identity

$$\int_{\mathbf{R}^N} A\mathbf{u} \cdot \mathbf{u} \, dx = \lambda \int_{\mathbf{R}^N} \left[ NF(x, \mathbf{u}) - \frac{1}{2} (N-2) \mathbf{u} \cdot \mathbf{f}(x, \mathbf{u}) + x \cdot \nabla_x F(x, \mathbf{u}) \right] dx. \quad (41)$$

The notation  $\nabla_x$  means that the components of  $\mathbf{u}$  are held constant during the differentiation. If  $a^i \equiv 0$  for each  $i = 1, \dots, M$ , then (41) holds for every solution  $\mathbf{u} \in D$  of (34). Similar Pohožaev identities are used in [6, 11, 13, 26, 30, 36], for example.

If  $\mathbf{u} \in E$  is a positive solution of (1) with the prototype potential (3), i.e.,

$$f^i(x, \mathbf{u}) = \sum_{k=1}^K \alpha_k^i g^k(x) (u^1)^{\alpha_k^1} \dots (u^i)^{\alpha_k^i - 1} \dots (u^M)^{\alpha_k^M},$$



then (41) implies that

$$0 \leq \int_{\mathbf{R}^N} \sum_{k=1}^K \left[ \left( N - \frac{1}{2}(N-2) |\alpha_k| \right) g^k(x) + x \cdot \nabla g^k(x) \right] u^{\alpha_k} dx. \quad (42)$$

If  $g^k(x) > 0$  and  $0 \neq x \cdot \nabla g^k(x) \leq 0$  throughout  $\mathbf{R}^N$  for  $k = 1, \dots, K$ , then necessarily  $|\alpha_k| < 2N/(N-2)$  for at least one integer  $k$ .

If  $|\alpha_k| = 2N/(N-2)$  for all  $k = 1, \dots, K$ , it also follows from (42) that  $x \cdot \nabla g^k(x)$  must be positive on a set of sufficiently large measure, for some  $k$ , in order that a positive  $\mathbf{u} \in E$  (or  $\mathbf{u} \in D$  if each  $a^i \equiv 0$ ) of (1) can exist.

We now derive necessary conditions for the existence of a positive solution  $\mathbf{u} \in D(\Omega)$  of (34) in an exterior domain  $\Omega$ . The conditions on  $\mathbf{f}$  stated below (34) will be retained, and we also need *lower* bounds

$$f^i(x, \mathbf{u}) \geq g_*^i(|x|) \mathbf{u}^{\beta_i}, \quad x \in \Omega, \quad 0 < \mathbf{u} \leq 1, \quad i = 1, \dots, M \quad (43)$$

for positive continuous functions  $g_*^i: \mathbf{R}_+ \rightarrow \mathbf{R}_+$ .

**THEOREM 10.** *Suppose  $f^i$  satisfies conditions  $(A_3)$  and (43), and  $\beta_i > 1$ . Then necessary conditions for the existence of a positive solution  $\mathbf{u} \in D(\Omega)$  of the system (34) are*

$$\int_0^\infty g_*^i(r) r^{N-1-\beta_i(N-2)} dr < \infty, \quad i = 1, \dots, M. \quad (44)$$

*Proof.* By Theorem 8,  $u^i(x) \rightarrow 0$  as  $|x| \rightarrow \infty$ , and hence the maximum principle for  $-\Delta u^i \geq 0$  implies that there exist positive constants  $C$  and  $R_0$  such that  $u^i(x) \geq C|x|^{2-N}$  for all  $|x| > R_0$ ,  $i = 1, \dots, M$ . It then follows from (34) and (43) that

$$-\Delta u^i \geq C^{M-1} p^i(|x|) (u^i)^{\beta_i}, \quad |x| \geq R_0, \quad i = 1, \dots, M, \quad (45)$$

where

$$p^i(r) = g_*^i(r) r^{-(\beta_i - 1)(N-2)}.$$

Since  $u^i > 0$ , a well known oscillation criterion for (45) [29] implies that

$$\int_0^\infty p^i(r) r^{N-1-\beta_i(N-2)} dr < \infty,$$

which is equivalent to (44). (This oscillation criterion has been rediscovered by many authors).

In particular, if (43) is specialized to

$$f^i(x, \mathbf{u}) \geq C^i |x|^{-b_i} \mathbf{u}^{\beta_i}, \quad x \in \Omega, \quad \mathbf{u} > 0, \quad i = 1, \dots, M,$$

for some constants  $b_i$  and  $C^i > 0$ , Theorem 10 shows that no positive solution  $\mathbf{u} \in D(\Omega)$  of (34) can exist if  $|\beta_i| \leq (N - b_i)/(N - 2)$  for some integer  $i$ .

The proof of Theorem 10 also establishes, without condition  $(A_3)$ , that (44) is a necessary condition for the existence of a minimal positive solution of (34) in  $\Omega$ , i.e., a solution  $\mathbf{u}$  such that  $|x|^{N-2} \mathbf{u}(x)$  is bounded above and below by positive constants in  $\Omega$ .

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